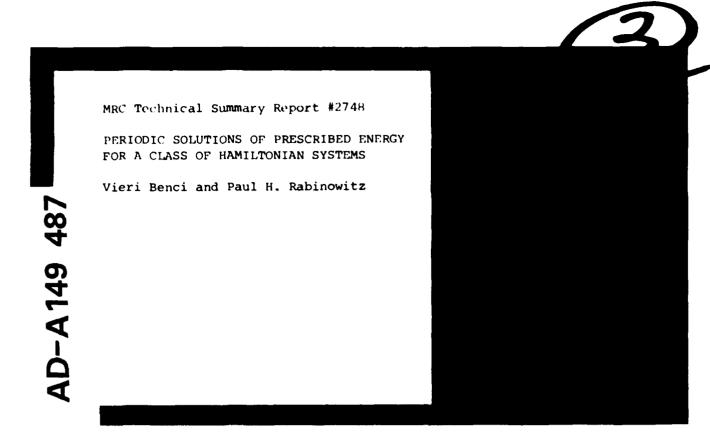


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PERIODIC SOLUTIONS OF PRESCRIBED ENERGY FOR A CLASS OF HAMILTONIAN SYSTEMS

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ABSTRACT

The main result of this paper is the following theorem: Let $p,q \in \mathbb{R}^n$, $H = H(p,q) \in C^2(\mathbb{R}^{2n},\mathbb{R})$ and let $H^{-1}(1)$ be the boundary of a compact neighborhood of 0 with $^{7}H \neq 0$ on $H^{-1}(1)$. If further $p \cdot H_p > 0$ on $H^{-1}(1)$ when $p \neq 0$, then the Hamiltonian system of ordinary differential equations

$$p = -H_q(p,q), \quad q = H_p(p,q)$$

possesses a periodic solution on $H^{-1}(1)$. The proof involves minimax arguments from the calculus of variations.

AMS (MOS) Subject Classifications: 34C15, 34C25

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SIGNIFICANCE AND EXPLANATION

Hamiltonian systems of ordinary differential equations model the motion of a discrete mechanical system when no frictional forces are present. A basic property of such systems is that "energy" is conserved. Therefore solutions of Hamiltonian systems lie on surfaces of fixed energy. The main result of this paper is a fairly general criterion for such a surface to possess a periodic solution.

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Introduction

Let $p,q \in \mathbb{R}^n$ and $H = H(p,q) : \mathbb{R}^{2n} + \mathbb{R}$ be smooth. The problem to be studied here is the existence of periodic solutions of the associated Hamiltonian system of ordinary differential equations

(HS)
$$\dot{p} = -\frac{\partial H}{\partial q}(p,q), \quad \dot{q} = \frac{\partial H}{\partial p}(p,q)$$

where $\cdot \equiv \frac{d}{dt}$. Setting z = (p,q), (HS) can also be written more succinctly as $z = JH_{x}(z)$

Theorem 1: Suppose H satisfies

(H₁) Hec²(\mathbb{R}^{2n} , \mathbb{R}),

(H₂) $\mathcal D$ is the boundary of a compact neighborhood of 0 and H₂ \neq 0 on $\mathcal D$ (i.e. $\mathcal D$ is a manifold).

 (H_3) p • $H_D \neq 0$ if $p \neq 0$.

Then (HS) possesses a periodic solution on \mathcal{D} .

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Theorem 1 has several predecessors. Seifert [1] considered Hamiltonians of the form

$$H(p,q) = \sum_{i,j=1}^{n} a_{ij}(q)p_{i}p_{j} + V(q)$$
,

i.e. H is a sum of kinetic and potential energy terms where $\mathcal{B} = \{q \in \mathbb{R}^n | V(q) \le 1\}$ is differentiable to the closed unit ball in \mathbb{R}^n and $\partial \mathcal{B}$ is a manifold, the matrix $\{a_1, a_2, a_3, a_4, a_5\}$ is uniformly positive definite in \mathcal{B} , and H is smooth. Using geodesic is unments from secondary, Seifert proved there exists a periodic solution of (HS) of a special type on \mathcal{B} . Generalizing his arguments, Weinstein [2] permitted a more general bicario energy term K(p,q) where for fixed q, K is even and convex in p while Gluck and Ziller [2] relaxed the condition on \mathcal{B} merely requiring \mathcal{B} to be compact with its poundary a manifold. See also Hayashi [4] and Benci [5] for results related to [3]. Another approach was made to (HS) in Rabinowitz [6] for H = K + V where V satisfied Siefert's condition and $p \in K_p > 0$ for $p \neq 0$. A case not covered by Theorem 1 but which can be obtained by similar but simpler arguments was given in [7] in which \mathcal{D} is the boundary of a compact star-shaped neighborhood of 0.

In a different direction from Theorem 1, there has been some recent work on the calciplicity of solutions of (HS) on \mathcal{D} , generally when \mathcal{D} bounds a convex region in \mathbb{R}^{2n} . See e.g. Ekeland-Lasry [8], Ambrosetti-Mancini [9], van Groesen [10], Berestycki-Lasry-Mancini-Ruf [11], and Ekeland [12].

We will prove Theorem 1 by a direct variational approach using minimax arguments. The proof relies in part on ideas from [5-6]. Let z(t) = (p(t),q(t)) be 2π periodic and

$$A(z) = \int_{0}^{2\pi} p \cdot qdt.$$

 $x = x^2 + y^2 +$

$$\Psi(z) = \frac{1}{2\pi} \int_{0}^{2\pi} \frac{1}{H(z)dt},$$

H(z) = H(z) on D and is suitably modified on $\mathbb{R}^{2n} \setminus D$. This critical point is produced as a minimax of $A|_{M}$ over an appropriate class of subsets of M. In this approach, the unknown period appears as a Lagrange multiplier.

The modified Hamiltonian H will be defined in §1 where some simple corollaries of Theorem 1 will also be obtained. In §2 the functional analytical framework in which the problem is treated is introduced. The properties of M and $A|_{M}$ such as the Palais-Smale condition are dealt with in §3. Theorem 1 is proved in §4. A dual variational argument is used in §5 to give an alternate approach to Theorem 1. In §6 a priori bounds from above and below are obtained for the unknown period of any solution of (HS) in terms of A(z). Lastly in §7, the results of §4 and 6 are used to prove a stronger version of Theorem 1 with (H₁) replaced by

(H)
$$H \in C^{1}(\mathbb{R}^{2n}, \mathbb{R})$$
.

An intriguing open question concerning (HS) is whether Theorem 1 remains true or is false if hypothesis (H_3) is omitted.

§1. The Modified Hamiltonian

For technical reasons that will become clear later, the Hamiltonian will be redefined outside of a neighborhood of \mathcal{D} . Suppose H satisfies $(H_1)^-(H_3)$. Then H(0) < 1 and H > 1 outside of the neighborhood of 0 bounded by \mathcal{D} . Without loss of generality we can assume $H(0) < \frac{1}{2}$. Our initial modification of H will allow us to assume H > 0, is a multiple of $|z|^2$ near 0, satisfies $(H_1)^-(H_3)$, and H_{ZZ} is uniformly bounded.

Indeed since $H(0) < \frac{1}{2}$, $\rho > 0$ can be chosen so that $\rho |z|^2 > \frac{3}{2}$ if H(z) > 1 and $\rho^{-1}|z|^2 < \frac{1}{4}$ if H(z) < 1. Let $\chi \in C^{\infty}(\mathbf{R},\mathbf{R})$ such that $\chi(s) = 0$ if $s < \frac{1}{4}$; $\chi(s) = 1$ if $s > \frac{1}{2}$, and $\chi'(s) > 0$ if $s \in (\frac{1}{4}, \frac{1}{2})$. Define

 $\widetilde{H}(z) = \chi(H(z) - 1)[\rho|z|^2 - H(z)] + H(z) + \chi(1 - H(z))[\rho^{-1}|z|^2 - H(z)].$ Then $\widetilde{H} \in C^2(\mathbb{R}^{2n}, \mathbb{R})$, $\widetilde{H} = \rho^{-1}|z|^2$ near z = 0, and $\widetilde{H} = H$ near \mathcal{D} . Moreover if $\widetilde{H}(z) = 1$, then $z \in \mathcal{D}$. To see this, suppose H(z) > 1. Then

$$\widetilde{H}(z) = \chi(H(z) - 1)[\rho|z|^2 - H(z)] + H(z)$$
.

If $H(z) > \frac{3}{2}$, $\widetilde{H}(z) = \rho |z|^2 > \frac{3}{2} > 1$ while if $H(z) \in (1, \frac{3}{2})$, $\widetilde{H}(z) > H(z) > 1$. Similar reasoning shows $\widetilde{H}(z) < 1$ if H(z) < 1. Thus $\widetilde{H}^{-1}(1) = \mathcal{D}$. A related argument shows $\widetilde{H}(z) > 0$ if $z \neq 0$. It is clear that \widetilde{H} satisfies $(H_1) - (H_2)$. To verify (H_3) , by the definition of χ , it suffices to show that for $p \neq 0$,

$$(1.1) \quad p \cdot \widetilde{H}_{p}(z) = \{\chi'(H(z) - 1)[\rho|z|^{2} - H(z)] - \chi'(1 - H(z))[\rho^{-1}|z|^{2} - H(z)]\}p \cdot H_{p}(z)$$

$$+ (1 - \chi(H(z) - 1) - \chi(1 - H(z)))p \cdot H_{p}(z)$$

$$+ 2|p|^{2}[\chi(H(z) - 1)\rho + \chi(1 - H(z))\rho^{-1}] > 0$$

if $z \in H^{-1}[\frac{1}{2}, \frac{3}{2}]$. Again this follows from our choice of ρ and (H_3) .

Since D is compact, there is a $\beta>0$ such that $|z|\leq\beta$ for $z\in D$. Let $\hat{\chi}\in C^\infty(\mathbf{R},\mathbf{R})$ such that $\hat{\chi}(s)=1$ for $s\leq 2\beta$, $\hat{\chi}(s)=0$ for $s\geq 4\beta$, and $\hat{\chi}'(s)<0$ for $s\in (2\beta,4\beta)$. Set

$$\hat{H}(z) = \hat{\chi}(|z|)\hat{H}(z) + (1 - \hat{\chi}(|z|))\hat{\rho}|z|^2$$
.

Then it is easy to check that for $\hat{\rho}$ chosen so that $\hat{\rho}|z|^2 > \hat{H}(z)$ for $|z| \in (2\beta, 4\beta)$, \hat{H}

possesses the properties verified above for \widetilde{H} , \widehat{H} is a multiple of $|z|^2$ for large |z|, and satisfies

 (H_4) $H_{zz}(z)$ is uniformly bounded.

Remark 1.2: The above arguments work equally well if H merely satisfies (H_1) and (H_3) in a neighborhood of \mathcal{D} .

Next hypothesis (H_3) will be used to decompose H into a sum of kinetic and potential energy terms. Set $U(q) = \hat{H}(0,q)$ and $K(p,q) \equiv K(z) \equiv \hat{H}(z) - U(q)$. Note that K(z) > 0 via (H_3) and $K, U \in \mathbb{C}^2$ via (H_3) . Moreover

Proposition 1.3: K satisfies the following properties:

$$(K_1) \quad K(0,q) = 0$$

$$(K_2)$$
 p • $K_p(z) > 0$ if $p \neq 0$

$$(K_3)$$
 $|K_D(z)| \le a_1(1 + |z|)$

$$(K_4)$$
 $K(z) \le a_1(1 + |z|)|p|$

$$(K_5)$$
 $|K_q(z)| \leq a_2|p|$

 (K_5) $K_{22}(z)$ is uniformly bounded in \mathbb{R}^{2n} .

(In (K_3) - (K_5) and later, a_i denotes a constant.)

<u>Proof:</u> (K_1) and (K_2) follow from the definition of K and (K_3) , and (K_3) , (K_6) from (H_4) . Since

(1.4)
$$K(p,q) = \int_{0}^{1} \frac{d}{ds} K(sp,q) ds = \int_{0}^{1} p \cdot K_{p}(sp,q) ds ,$$

 (K_3) and (1.4) imply (K_4) . Similarly

(1.5)
$$K_{q}(p,q) = H_{q}(p,q) - H_{q}(0,q) = \int_{0}^{1} H_{pq}(sp,q)p \ ds$$

so (H_4) and (1.5) give (K_5) .

To define H, one final modification of H is required. Let

$$\Omega_{\mathbf{q}} = \{ \mathbf{q} \in \mathbb{R}^n | \mathbf{U}(\mathbf{q}) < \mathbf{s} \}$$
.

By (H₂), $U_q \neq 0$ on $\partial \Omega_1$ and there exist constants $d, \beta > 0$ such that $U_q(q) \neq 0$ and (1.6) $|U_q(q)| > \beta U(q) \text{ if } q \in \Omega_{1+2d} \setminus \Omega_{1+2d} .$

Let $\phi \in \mathbb{C}^2$ be defined for s < 1 + 2d such that

$$(\phi_1)$$
 $\phi(s) = s$, $s \le 1 + d$

$$(\phi_2)$$
 $\phi^1(s) > 1$, $s < 1 + 2d$

$$(\phi_3)$$
 $\phi(s) = (s - (1 + 2d))^{-2}$, s near 1 + 2d.

te further extend ϕ to all of R via $\phi(s) = \infty$ if s > 1 + 2d. Finally define $V(q) \equiv \phi(U(q))$ for $q \in \mathbb{R}^n$ and

$$H(z) \equiv K(z) + V(q)$$
.

Thus \overline{H} is C^2 where finite and if $\overline{H}(z) = 1$, $U(\sigma) \leq V(q) \leq 1$ by (ϕ_2) . Thus V(q) = U(q) by (ϕ_1) and $\overline{H}(z) = H(z)$. Consequently $\overline{H}^{-1}(1) = \mathcal{D}$.

We will find a periodic solution of

(1.7)
$$p = -H_q$$
, $q = H_p$

on $H^{-1}(1)$. Hence it will be a periodic solution of (HS) on \mathcal{D} .

To conclude this section, some estimates will be obtained for V. Let $\psi \in C^{\infty}(R,R)$ such that

$$\psi(s) = 0, \quad s < 1 - 2d$$

and $\psi'(s) > 0$ if $s \in (1 - 2d, 1 - d)$. Set

$$v(q) = \psi(v(q)) \frac{u_{q}(q)}{|u_{q}(q)|}.$$

Observing that the ψ term vanishes if $q \in \Omega_{:-2d}$ and $\psi'(V(q))$ vanishes if V(q) > 1 - d, it follows that $V \in C^1(\mathbb{R}^n, \mathbb{R})$.

Proposition 1.8: There is a constant $\gamma > 0$ such that

(1.9)
$$v_{\alpha}(q) \cdot v(q) > 0$$
 for all $q \in \Omega_{1+2d}$

and

(1.10)
$$V_q(q) \cdot V(q) > \gamma V(q)$$
 for all $q \in \Omega_{1+2d} \setminus \Omega_{1-2d}$

<u>Proof:</u> Inequality (1.9) is immediate from the definition of V and (ϕ_2) . To check (1.10), note that if $U(q) \in (1-d,1+2d)$, then V(q) > U(q) by (ϕ_2) and $\psi(V(q)) = 1$. Therefore

(1.11)
$$v_{q}(q) \cdot v(q) = \phi'(v(q))|v_{q}| > \beta(1 - d)\phi'(v(q))$$

by (1.6). Thus to get (1.10), it suffices to show

(1.12)
$$\phi'(s) \ge \beta_1 \phi(s)$$
, $s \in (1 - d, 1 + 2d)$

and this is immediate from (ϕ_1) - (ϕ_3) .

§2. Functional Analytical Preliminaries

The space in which (1.7) will be treated is the Hilbert space

$$E = \{z = (p,q) | p \in L^2(S^1,R^N), q \in W^{1,2}(S^1,R^N)\}$$

 $= L^2(S^1,R^N) \oplus W^{1,2}(S^1,R^N)$

where $L^2(S^1,\mathbb{R}^n)$ denotes the set of n-tuples of 2π periodic functions which are square integrable, etc. For we $L^2(S^1,\mathbb{R}^n)$, let

$$[w] = \frac{1}{2\pi} \int_{0}^{2\pi} w(t) dt.$$

Thus any $z = (p,q) \in E$ can be decomposed into $(\{p\},[q]) + (\hat{p},\hat{q})$ where $(\hat{p},\hat{q}) \in \hat{L}^2 \oplus \hat{w}^{1,2}$ and

$$\hat{L}^{2} = \{ p \in L^{2}(S^{1}, R^{n}) | \{ p \} = 0 \} ,$$

$$\hat{w}^{1,2} = \{ q \in W^{1,2}(S^{1}, R^{n}) | \{ q \} = 0 \} .$$

As inner product in E we take

$$(z_1, z_2)_E = \int_0^{2\pi} [(\hat{p}_1(t) \cdot \hat{p}_2(t)) + (\hat{p}_1 \cdot \hat{p}_2)]dt + [p_1][p_2] + [q_1][q_2]$$

where $D = \frac{d}{dt}$ and $z_1 = (p_1, q_1) = ([p_1] + \hat{p}_1, [q_1] + \hat{q}_1)$, etc. The norm in E will be denoted by $\| \cdot \|$ and we will generally use the same notation for the norm in E^* , the dual space of E.

It is easy to see that $D|_{\widehat{W}_{1,2}}:\widehat{W}^{1,2}+\widehat{L}^2$ is an isomorphism. Let D^{-1} denote its inverse. We define linear maps P^0,P^+,P^- of E into E by

$$P^{0}(p,q) \equiv ([p],[q])$$

and

$$p^{\pm}(p,q) \equiv (\frac{1}{2} (\hat{p} \pm p\hat{q}), \frac{1}{2} (\hat{q} \pm p^{-1}\hat{p}))$$
.

It is easy to verify that these maps are well defined and are (continuous) projectors on

E satisfying $p^0 + p^+ + p^- = id$, the identity map on E. Define $E^0 \equiv p^0 E$ and $E^{\pm} \equiv P^{\pm} E$. Note that if $(p^{\pm}, q^{\pm}) \in E^{\pm}$, then $p^{\pm} = \frac{1}{2} (p^{\pm} \pm Dq^{\pm})$. Therefore (2.1)

Next observe that the spaces E^0 , E^{\pm} are mutually orthogonal subspaces of E. E.g. if $z^{\pm}=(p^{\pm},q^{\pm})$ \in E^{\pm} ,

$$(z^{+},z^{-})_{E} = \int_{0}^{2\pi} [(p^{+} \cdot p^{-}) + (Dq^{+} \cdot Dq^{-})]dt$$

$$= \int_{0}^{2\pi} [(Dq^{+} \cdot (-Dq^{-})) + (Dq^{+} \cdot Dq^{-})]dt = 0$$

via (2.1).

For $z = (p,q) \in E$, define the action integral as

$$A(z) \equiv \int_{0}^{2\pi} p \cdot \dot{q} dt.$$

Then $A \in C^{\infty}(E,R)$ and writing $z = z^0 + z^+ + z^-$ and using (2.1) shows

2.2)
$$A(z) = \int_{0}^{2\pi} (p^{0} + p^{+} + p^{-}) \cdot (pq^{+} + pq^{-}) dt$$

$$= \int_{0}^{2\pi} \{(p^{+} \cdot pq^{+}) + (p^{+} \cdot pq^{-}) + (p^{-} \cdot pq^{+}) + (p^{-} \cdot pq^{-})\} dt$$

$$= \frac{1}{2} \int_{0}^{2\pi} (|p^{+}|^{2} + |pq^{+}|^{2} - |p^{-}|^{2} - |pq^{-}|^{2}) dt$$

$$= \frac{1}{2} (||z^{+}||^{2} - ||z^{-}||^{2}) .$$

Next define

$$\Psi(z) \equiv \frac{1}{2\pi} \int_{0}^{2\pi} H(z) dt$$

and

$$M \equiv \{z \in E | \Psi(z) = 1\}.$$

Our goal is to obtain a periodic solution of (HS) (or equivalently (1.7)) as a critical point of $A|_{M}$. As will be seen later, a critical point z of this constrained variational problem satisfies z e $C^1(S^1,R^{2n})$ and

(2.3)
$$\dot{z} = \lambda J H_z(z)$$

where $\lambda \in \mathbb{R} \setminus \{0\}$. Since (2.3) is a Hamiltonian system, $H(z(t)) \equiv \text{constant}$. Thus $z \in M$ implies $z \in \mathcal{D}$. Moreover since $\lambda \neq 0$, rescaling time in (2.3) yields a periodic solution of (1.7) on \mathcal{D} , i.e. the desired solution of (HS).

§3. Some Properties of M and A

This section studies several properties of Ψ and M. It will be shown that M is a $C^{1,1}$ manifold which bounds a neighborhood of 0 in E and $A|_{M}$ satisfies a version of the Palais-Smale condition.

For $q \in W^{1,2}(S^1, \mathbb{R}^n)$, set

$$V(q) = \frac{1}{2\pi} \int_{0}^{2\pi} V(q(t))dt$$
.

Remark 3.1: The definition of V shows that if $V(q) < \infty$, $q(t) \in \Omega_{1+2d}$ for almost all $t \in [0,2\pi]$ and since $q \in C(S^1,\mathbb{R}^n)$, $q(t) \in \overline{\Omega}_{1+2d}$ for all $t \in [0,2\pi]$. In particular there is a constant M > 0 (and independent of q) such that $\|q\|_{\infty} < M$.

For $x \in \mathbb{R}^n$, let

$$\ell(x) \equiv \inf_{y \in \partial \Omega_{1+2d}} |x - y|.$$

Proposition 3.2: There exist constants β_1, M_1 such that for all $x \in \overline{\Omega}_{1+2d}$ (3.3) $\beta_1 \ell(x)^{-2} \leq V(x) + M_1$

<u>Proof:</u> Let $x \in \Omega_{1+2d}$. Using the implicit function theorem, it is not difficult to show that there is an $\epsilon_0 > 0$ such that if $\ell(x) < \epsilon_0$, there exists a unique $x \in \partial \Omega_{1+2d}$ and $\rho > 0$ such that

$$x = x - \rho u_q(\bar{x})$$
.

Therefore there is a $\beta_2 > 0$ such that

(3.4)
$$\beta_2 \rho > |x - \overline{x}| > \beta \rho \text{ if } \ell(x) < \epsilon_0$$

via the continuity of $U_{\mathbf{q}}$ and (1.6). Now

(3.5)
$$U(x) = \overline{U(x)} = \overline{U_q(x)}(x - x) + o(|x - x|)$$

as $x + \partial \Omega_{1+2d}$. Therefore by (3.4) - (3.5), for x near $\partial \Omega_{1+2d}$, e.g. $\ell(x) < \epsilon$,

(3.6)
$$|U(x) - U(\bar{x})| \le \rho |U_q(\bar{x})|^2 + o(q) \le M_2 \rho$$
.

Now for $l(x) \le \varepsilon$, by (ϕ_3) and (3.6),

$$V(x) = (U(x) - (1 + 2d))^{-2} > (M_2 \rho)^{-2}$$
.

But $\rho = \ell(x) |U_{\sigma}(\bar{x})|^{-1}$. Therefore

$$V(x) > M_3 \ell(x)^{-2}$$

if $l(x) \le \varepsilon$. If $l(x) > \varepsilon$, $l(x)^{-2} \le \varepsilon^{-2}$ so

$$v(x) + M_3 \varepsilon^{-2} > M_3 \ell(x)^{-2}$$
.

Thus (3.3) obtains with $M_2 = M_3 \epsilon^{-2}$ and $\beta_1 = M_3$.

The estimate (3.3) will be used next to show that $V(q) < \infty$ implies that q(t) avoids $\partial \Omega_{1+2d}$.

Proposition 3.8: Let $q \in W^{1,2}(S^1, \mathbb{R}^n)$ satisfy $V(q) < \infty$. Then there is an $\widetilde{M} = \widetilde{M}(\|q\|_{W^{1,2}}, V(q)) > 0$ such that $\ell(q(t)) > \widetilde{M}$ for all $t \in [0, 2\pi]$.

Proof: Since q is 2m periodic, by translating t it can be assumed that

$$\ell(q(0)) = \min_{t \in [0,2\pi]} \ell(q(t)) \equiv \mu.$$

By the Cauchy-Schwarz inequality,

(3.9)
$$|q(t) - q(0)| \le \int_{0}^{t} |\dot{q}(\tau)| d\tau \le t^{1/2} ||q||_{\dot{W}^{1,2}}.$$

Since & is Lipschitz continuous (with constant 1),

(3.10)
$$|\ell(q(t)) - \ell(q(0))| \le |q(t) - q(0)| \le t^{1/2} ||q||_{t_0^{1/2}}.$$

Therefore

(3.11)
$$\ell(q(t)) \leq \mu + t^{1/2} |q|_{W^{1,2}} .$$

We can assume $\|q\|_{W^{1,2}} > 0$ for otherwise the result is trivial. By (3.3) and (3.11),

(3.12)
$$\Rightarrow \frac{1}{2\pi} \int_{0}^{2\pi} (\beta_{1} \ell(q(t))^{-2} - M_{1}) dt$$

$$\Rightarrow \frac{\beta_{1}}{2\pi} \int_{0}^{2\pi} (\mu + t^{1/2} \|q\|_{W^{1}, 2})^{-2} dt - M_{1}$$

$$\Rightarrow \frac{\beta_{1}}{\pi} \int_{0}^{2\pi} (\mu^{2} + t \|q\|_{W^{1}, 2}^{2})^{-1} dt - M_{1}$$

$$= \frac{\beta_{1}}{\pi \|q\|_{W^{1}, 2}^{2}} \log(1 + \frac{2\pi \|q\|_{W^{1}, 2}^{2}}{\mu^{2}}) - M_{1}$$

from which the result follows.

Remark 3.13: Proposition 3.8 implies that the domain of V is $\{q \in W^{1,2}(S^1,\mathbb{R}^n) | q(t) \in \Omega_{1+2d} \text{ for all } t \in [0,2\pi] \}.$

The smoothness of V will be established next.

Proposition 3.14: $V \in \mathbb{C}^2$ on the domain of V.

Proof: Let $q \in W^{1,2}(S^1,\mathbb{R}^n)$ with $V(q) < \infty$. Let $\delta = \inf_{\substack{t \in [0,2\pi] \\ W^{1,2}}} \ell(q(t))$. Then $\delta > 0$ by $\ell(q) = \inf_{\substack{t \in [0,2\pi] \\ W^{1,2}}} \ell(q)$. Therefore since $V \in C^2(\Omega_{1+2d},\mathbb{R})$,

$$(3.15) \qquad \qquad V(q+\widetilde{q}) = V(q) + V_q(q)\widetilde{q} + \frac{1}{2} V_{qq}(q)(\widetilde{q},\widetilde{q}) + o(|\widetilde{q}|^2)$$
 as $\widetilde{q} \neq 0$ uniformly for te $[0,2\pi]$. The definition of V , (3.15) , and the compact embedding of $W^{1,2}(S^1,\mathbb{R}^n)$ in $C(S^1,\mathbb{R}^n)$ then readily imply that V is Frechet differentiable at q with

$$V'(q)\tilde{q} = \frac{1}{2\pi} \int_{0}^{2\pi} V'(q)\tilde{q} dt$$
,

V'(q) is continuous, $(V'(q))' \equiv V''(q)$ exists,

$$V^{*}(q)(\tilde{q},\tilde{q}) = \frac{1}{2\pi} \int_{0}^{2\pi} V^{*}(q)(\tilde{q},\tilde{q})dt$$

and $V^{n}(q)$ is continuous.

For z e E, set

$$K(z) = \frac{1}{2\pi} \int_{0}^{2\pi} K(z) dt .$$

Proposition 3.16: $K \in \mathbb{C}^{1,1}(E,\mathbb{R})$, (i.e. K is Frechet differentiable and its Frechet derivative is Lipschitz continuous).

Proof: Since $K \in C^2(\mathbb{R}^{2n}, \mathbb{R})$, given any $z, \zeta \in \mathbb{R}^{2n}$, by Taylor's Theorem,

(3.17)
$$K(z + \zeta) = K(z) + K_{z}(z)\zeta \frac{1}{2}K_{zz}(z + \theta\zeta)(\zeta,\zeta)$$

for some $\theta \in (0,1)$. By (K_6) of Proposition 1.3, K_{ZZ} is uniformly bounded. Therefore there is a constant $M_3 > 0$ such that

$$|\kappa(z+\zeta)-\kappa(z)-\kappa_z(z)\zeta| \leq M_3|\zeta|^2$$

for all $z,\zeta \in \mathbb{R}^{2n}$. Choosing $z,\zeta \in E$, (3.18) implies

(3.19)
$$|K(z + \zeta) - K(z)| - \frac{1}{2\pi} \int_{0}^{2\pi} K_{z}(z) \zeta \, dt | \leq \frac{M_{3}}{2\pi} ||\zeta||_{L^{2}}^{2} \leq M_{4} ||\zeta||^{2}.$$

In particular for z fixed, given any $\varepsilon > 0$, if ζ is sufficiently small, the right hand side of (3.19) does not exceed $\varepsilon \| \zeta \|$. Hence K is Frechet differentiable and

$$K'(z)\zeta = \frac{1}{2\pi} \int_{0}^{2\pi} K_{z}(z)\zeta dt.$$

To show that K' is Lipschitz continuous, note that

(3.20)
$$|K'(z+w) - K'(z)| = \sup_{z \in \mathcal{L}, |z| \le 1} \left| \frac{1}{2\pi} \int_{0}^{2\pi} (K_{z}(z+w) - K_{z}(z)) \zeta dt \right| .$$

As in (3.17) by the Mean Value Theorem and (K_6) ,

(3.21)
$$|K_{z}(z+w) - K_{z}(z)| \le M_{5}|w|$$

for some constant M_5 . Therefore (3.20)-(3.21) imply

(3.22)
$$||K'(z + w) - K'(z)||_{\Phi} \leq M_6 ||w||,$$

i.e. K' is Lipschitz continuous.

By the definitions of K and V, $\Psi = K + V$. Recall that $M = \Psi^{-1}(1)$. The next three propositions study some properties of M.

Proposition 3.23: M is C1,1 manifold in E.

<u>Proof:</u> The smoothness assertions follow on combining Proposition 3.14 and 3.16 once we show that M is a manifold, i.e. $\Psi^*(z) \neq 0$ for all $z \in M$. But if z = (p,q) with $p \not\equiv 0$,

$$\Psi^{*}(z)(p,0) = \frac{1}{2\pi} \int_{0}^{2\pi} p \cdot K_{p}(z)dt > 0$$

via (K_2) of Proposition 1.3. If $p\equiv 0$, then $\Psi(z)=V(q)=1$ via (K_1) and by Proposition 3.8, $q(t)\in\Omega_{1+2d}$ for all $t\in [0,2\pi]$. If $q(t)\in\Omega_{1-d}$ for all $t\in [0,2\pi]$, $V(q)\leq 1-d$ which is impossible. Thus since $q\in C(S^1,\mathbb{R}^n)$, $q(t)\in\Omega_{1+2d}\backslash\Omega_{1-d}$ on a set Y of positive measure. By previous remarks, $V(q)\in C(S^1,\mathbb{R}^n)$ and $(0,V(q))\in E$. Hence by (1.9)-(1.10), $\Psi^*(z)(0,V(q))={}^*(q)V(q)\geq \frac{1}{2\pi}\int_{\mathbb{R}^n} \gamma(1-d)dt>0$.

Thus M is a manifold and the proof is complete.

Proposition 3.25: M is the boundary of a neighborhood of 0 in E.

<u>Proof:</u> Ψ is continuous on $L^2(S^1,\mathbb{R}^n)$ e $\{q \in W^{1,2}(S^1,\mathbb{R}^n) | V(q) < \infty\}$ which is an open set in E. Therefore $\Psi^{-1}(-\infty,1)$ is open. Since H(0) = 0, 0 belongs to this set. Proposition 3.26: M is bounded in $L^2(S^1,\mathbb{R}^{2n})$.

<u>Proof:</u> Let $(p,q) \in M$. Since $V(q) < \infty$, by Remark 3.1, there is an M > 0 such that $\|q\|_{\infty} < M$. The definition of K implies it is a multiple of $\|p\|^2$ for $\|z\|$ near ∞ . L

$$K(z) > M_7 |_P|^2 - M_8$$

for all $z \in \mathbb{R}^{2n}$. Therefore

$$K(z) > \frac{M_7}{2\pi} ||p||_{L^2}^2 - \frac{M_8}{2\pi}.$$

Thus if zeM,

$$I_{P}I_{L^{2}}^{2} \le \left(1 + \frac{M_{B}}{2\pi}\right) \frac{2\pi}{M_{7}}$$
.

The next proposition shows that $A \mid_{M}$ satisfies a version of the Palais-Smale condition.

Proposition 3.27: A \mid_{M} satisfies (PS) $^{+}$, i.e. if c > 0 and (z_{j}) is a sequence in E such that

- (i) z; e M,
- (ii) $A(z_i) + c$,

and (iii) $A'(z_j) \sim \lambda_j \Psi'(z_j) \neq 0$ (in E^*)

as j → ∞ where

$$\lambda_{j} = (A'(z_{j}), \Psi'(z_{j}))_{E^{*}} \Psi'(z_{j}) \Psi^{-2}$$

then (z_j) has a convergent subsequence.

<u>Proof</u>: By (ii), there is an $\varepsilon \in (0, \frac{c}{2})$ such that

(3.28)
$$c - \varepsilon \leq A(z_j) \leq c + \varepsilon$$

for all large $j \in \mathbb{N}$. Similarly by (iii), there is a $w_j \in E^*$ with $w_j + 0$ as $j + \infty$ and

(3.29)
$$A'(z_{ij})\zeta - \lambda_{ij}\Psi'(z_{ij})\zeta = \langle w_{ij},\zeta \rangle$$

where $\langle \cdot, \cdot \rangle$ denotes the duality between E^* and E. Choosing $\zeta = (p_i, 0)$ yields

(3.30)
$$A(z_{j}) - \frac{\lambda_{j}}{2\pi} \int_{0}^{2\pi} \kappa_{p}(z_{j}) p_{j} dt = \langle w_{j}, (p_{j}, 0) \rangle.$$

By (3.28) and the choice of ϵ , $A(z_i) > \frac{c}{2}$ while the left hand side of (3.30) goes to 0

as $j + \infty$ since $((p_j, 0))$ is bounded in E. Consequently by $(K_2)-(K_3)$ of Proposition 1.3, λ_j is positive and bounded away from 0 for large j. We restrict ourselves to such $j \in \mathbb{N}$.

Choosing $\zeta = (\xi, \eta)$, (3.29) can be rewritten as

(3.31)
$$\int_{0}^{2\pi} [(p_{j} \cdot \hat{\eta}) + (\xi \cdot \hat{q}_{j})] dt - \frac{\lambda_{j}}{2\pi} \int_{0}^{2\pi} [(\tilde{H}_{p}(z_{j}) \cdot \xi) + (\tilde{H}_{q}(z_{j}) \cdot \eta)] dt = \langle w_{j}, \zeta \rangle.$$

Setting $w_j = (u_j, v_j) \in E^*$, i.e. $u_j \in L^2(S^1, \mathbb{R}^n)$ and $v_j \in W^{-1,2}(S^1, \mathbb{R}^n)$ and noting that $p_j \in W^{-1,2}(S^1, \mathbb{R}^n)$, (3.31) implies that

(3.32) (i)
$$-\dot{p}_{j} = \frac{\lambda_{j}}{2\pi} H_{q}(z_{j}) + v_{j}$$

(ii)
$$q_j = \frac{\lambda_j}{2\pi} H_p(z_j) + u_j = \frac{\lambda_j}{2\pi} K_p(z_j) + u_j$$

holds in the sense of distributions.

We claim that

(3.33)
$$\int_{0}^{2\pi} K_{\mathbf{p}}(\mathbf{z}_{\mathbf{j}}) \cdot \mathbf{p}_{\mathbf{j}} dt > a > 0$$

for all j \in N. Assuming (3.33) for now, (3.30) and (3.28) then imply λ_j is bounded away from ∞ . Then by (3.32) (ii),

so (K_3) of Proposition 1.3, the boundedness of (z_j) in L^2 (via Proposition 3.26) and Remark 3.1 show that (q_j) is bounded in $W^{1,2}(S^1,\mathbb{R}^n)$. Proposition 3.8 then implies there is an M>0 independent of j such that $L(q_j(t))>M$ for all $t\in[0,2\pi]$, i.e. the functions q_j lie uniformly inside Ω_{1+2d} . Therefore the functions $(V_q(q_j))$ are bounded in L^∞ and by (K_5) of Proposition 1.3, $H_q(z_j)$ are bounded in L^2 . Consequently

the right hand side of (3.32)(i) converges strongly in $W^{-1,2}(S^1,\mathbb{R}^n)$ (along a subsequence). Therefore p_j converges strongly in $W^{-1,2}(S^1,\mathbb{R}^n)$. Consequently $D^{-1}p_j$ converges strongly in $L^2(S^1,\mathbb{R}^n)$. But $p_j=\{p_j\}+D^{-1}p_j$. Hence along a subsequence p_j converges strongly in $L^2(S^1,\mathbb{R}^n)$. Lastly by (3.32) (ii), the same is true for a subsequence of q_j in $W^{1,2}(S^1,\mathbb{R}^n)$ since (K_3) of Proposition 1.3 implies that K_p is a continuous map of $L^2(S^1,\mathbb{R}^{2n})$ to $L^2(S^1,\mathbb{R}^n)$.

Thus Proposition 3.27 will be established once we show that (3.33) holds. Suppose that (3.33) is false, i.e.

(3.35)
$$\int_{0}^{2\pi} \kappa_{p}(z_{j}) \cdot p_{j}dt + 0$$

for some subsequence of j's. This implies that $\|\mathbf{p}_j\|_{L^2} + 0$ along this subsequence. Indeed let $\mathbf{Y}_{1j} = \{\mathbf{t} \in [0,2\pi] \, | \, |\mathbf{p}_j(\mathbf{t}) \, | < \sigma \}$, $\mathbf{Y}_{2j} = \{\mathbf{t} \in [0,2\pi] \, | \, |\sigma < |\mathbf{p}_j(\mathbf{t}) \, | < 4\beta \}$, and $\mathbf{Y}_{3j} = \{\mathbf{t} \in [0,2\pi] \, | \, |\mathbf{p}_j(\mathbf{t}) \, | > 4\beta \}$ where $\sigma < 1$ and β was defined together with $\hat{\mathbf{H}}$ in §1. On \mathbf{Y}_{3j} , $\mathbf{p}_j \cdot \mathbf{K}_p(\mathbf{z}_j) = 2\tilde{\rho} |\mathbf{p}_j|^2$. On \mathbf{Y}_{2j} , both $|\mathbf{p}|^2$ and $\mathbf{p} \cdot \mathbf{K}_p$ are bounded away from 0. Therefore there is a constant $\alpha = \alpha(\sigma)$ such that $|\mathbf{p}|^2 < \alpha(\sigma)\mathbf{p} \cdot \mathbf{K}_p$ on \mathbf{Y}_{2j} . Combining these observations yields

Since σ is arbitrary, (3.35)-(3.36) show $\|\mathbf{p}_j\|_{\mathbf{L}^2} \to 0$ as $j \to \infty$. Then (κ_3) - (κ_4) of Proposition 1.3 imply

(3.37)
$$K(z_1) + 0$$

as $j \rightarrow \infty$ while (K_3) and Proposition 3.26 show

(3.38)
$$\int_{0}^{2\pi} |\kappa_{p}(z_{j})| |p_{j}| dt + 0$$

as $j + \infty$. Also by (K_5) ,

(3.39)
$$|\int_{0}^{2\pi} K_{q}(z_{j}) \cdot \nu(q_{j}) dt| \leq |K_{q}(z_{j})|_{L^{2}} |\nu(q_{j})|_{L^{2}} + 0$$

as j $\rightarrow \infty$ since ($\|v(q_j)\|_{L^2}$) is uniformly bounded. By (3.37) and (i) of Proposition 3.27,

$$(3.40)$$
 $V(g_4) + 1$

as j + . Set

$$\Delta_{j} \equiv \{t \in [0,2\pi] | v(q_{j}(t)) > 1 - d \}$$
.

Therefore

(3.41)
$$V(q_{j}) \leq \frac{1}{2\pi} \left(\int_{\Delta_{j}} V(q_{j}(t)) dt + \int_{[0,2\pi] \setminus \Delta_{j}} V(q_{j}(t)) dt \right)$$

$$\leq \frac{1}{2\pi} \int_{\Delta_{j}} V(q_{j}(t)) dt + 1 - d.$$

Hence for large j, by (3.40)-(3.41),

$$\frac{\mathrm{d}}{2} < \frac{1}{2\pi} \int_{\Delta_{j}} v(q_{j}(t)) dt .$$

Next by (3.32) (1) and (1.9),

$$(3.43) \qquad \int_{0}^{2\pi} p_{j} \cdot \frac{d}{dt} \nu(q_{j}(t))dt = \frac{\lambda_{j}}{2\pi} \int_{0}^{2\pi} \overline{H}_{q}(z_{j}) \cdot \nu(q_{j})dt + \langle v_{j}, \nu(q_{j}) \rangle$$

$$> \frac{\lambda_{j}}{2\pi} \left[\int_{\Delta_{j}} v_{q}(q_{j}) \cdot \nu(q_{j})dt - \|K_{q}(z_{j})\|_{L^{2}} \|\nu(q_{j})\|_{L^{2}} \right] - \|v_{j}\|_{W^{-1,2}} \|\nu(q_{j})\|_{W^{1,2}}.$$

As has been noted earlier, $v(\cdot) \in C^1(\mathbb{R}^n,\mathbb{R})$ and $\|q\|_{\infty} \leq M$ for all $(p,q) \in M$. Hence there exists a constant $M_1 > 0$ and independent of j such that

$$||v(q_j)||_{W^{1,2}} \le M_1(1 + ||\dot{q}_j||_{L^2})$$
.

Thus (3.42)-(3.43), (1.10), and (K_5) of Proposition 1.3 imply that

as $j + \infty$. On the other hand, by (3.32) and (3.38),

$$\langle \lambda_{j} \circ (1) + \circ (1)^{2}$$

as $j \rightarrow \infty$. Combining (3.44)-(3.45) shows

(3.46)
$$\frac{\gamma d}{4} \lambda_{j} \leq o(1) \|\dot{q}_{j}\|_{L^{2}} + o(1)$$

as $j + \infty$. But then by (3.32) (ii),

(3.47)
$$\frac{\gamma d}{4} \lambda_{j} \leq o(1) \lambda_{j} ||K_{p}(z_{j})||_{L^{2}} + o(1)$$

so (K_3) of Proposition 1.3 and Proposition 3.26 imply $\lambda_j + 0$ as $j + \infty$, a contradiction. Thus (3.33) has been verified and Proposition 3.27 has been established.

Two further technical results are needed in this section. Let ℓ denote the duality map between $w^{-1,2}(s^1,\mathbb{R}^n)$ and $w^{1,2}(s^1,\mathbb{R}^n)$, i.e. ℓ is defined by

$$(Lw,\xi)_{w^{1,2}} = \langle w,\xi \rangle$$

for $w \in W^{-1,2}(S^1,\mathbb{R}^n)$ and $\xi \in W^{1,2}(S^1,\mathbb{R}^n)$. Abusing notation somewhat, we will also let L denote the duality between E and E^* . For $z = (p,q) \in E$, define $P_1z \equiv p$ and $P_2z \equiv q$.

Proposition 3.48: $P_2 L^{\Psi^*}$ is a compact map of M into $W^{1,2}(S^1,\mathbb{R}^n)$.

<u>Proof</u>: For $z \in M$ and $\zeta = (u,v) \in E$,

$$\Psi^{*}(z)\zeta \approx \langle \Psi^{*}(z),\zeta \rangle = \frac{1}{2\pi} \int_{0}^{2\pi} [(\vec{H}_{p}(z) \cdot u) + (\vec{H}_{q}(z) \cdot v)]dt = (\zeta \Psi^{*}(z),\zeta)_{E}.$$

Therefore

(3.49)
$$(P_2L\Psi^*(z),v)_{W^1,2} = \frac{1}{2\pi} \int_{0}^{2\pi} \bar{H}_{q}(z) \cdot vdt.$$

The right hand side of (3.49) is a continuous linear functional on $w^{1,2}(S^1,\mathbb{R}^n)$. Therefore there exists a unique $\theta = \theta(z) \in w^{1,2}(S^1,\mathbb{R}^n)$ such that

(3.50)
$$(\Theta(z), v)_{W_{1,2}} = \frac{1}{2\pi} \int_{0}^{2\pi} \tilde{H}_{q}(z) \cdot vdt .$$

Clearly $\Theta(z) = P_2 L\Psi^*(z)$. But the map $z \to H_q(z)$, $M \to L^2(S^1, \mathbb{R}^n)$ is continuous and the map $H_q(z) + \Theta(z)$, $L^2(S^1, \mathbb{R}^n) \to W^{1,2}(S^1, \mathbb{R}^n)$ is compact and linear. It follows that $P_2 L\Psi^*$ is compact.

The final result in this section is a version of the so-called "Deformation Theorem" which is appropriate for our setting. A subset $S \subset E$ will be called invariant if $z(t) \in S$ implies that $z(t+\theta) \in S$ for all $t \in [0,2\pi]$. A mapping h: S + E, where S is invariant, will be called equivariant if $h(T_{\theta}z) = T_{\theta}h(z)$ for all $\theta \in [0,2\pi]$ where $T_{\theta}\zeta(t) = \zeta(t+\theta)$. For $s \in R$, let

$$A_{s} \equiv \{z \in M | A(z) > s\} .$$

For ceR, let

$$K_{c} = \{z \in M | A(z) = c \text{ and } A'(z) = (A'(z), \Psi'(z)) + \|\Psi'(z)\|^{-2} \Psi'(z)\}$$

i.e. K_C is the set of critical points of $A|_M$ having critical value c.

Proposition 3.51: Let $c, \overline{c} > 0$. Then there is an $\varepsilon \in (0, \overline{c})$ and $\eta \in C([0,1] \times M, M)$ such that

1⁰ η(1,•) is equivariant

$$2^{\circ}$$
 $\eta(1,z) = z$ if $A(z) \notin [c - \overline{\epsilon}, c + \overline{\epsilon}]$

$$3^{\circ}$$
 $\ln(1,z) - z \leq 1$

$$4^{\circ}$$
 If $K_{c} = \emptyset$, $\eta(1, A_{c-\epsilon}) \subset A_{c+\epsilon}$

$$5^{\circ}$$
 $P^{\dagger} \eta(1,z) = \beta^{\dagger}(z)z^{\dagger} + B^{\dagger}(z)$ where $\beta^{\dagger} \in C(M,[1,e])$ and $P_{2}B^{\dagger}$ is compact.

<u>Proof:</u> Most of the above assertions follow from standard arguments and therefore we will be somewhat sketchy below. See e.g. [13-15] for more details. The function η is determined as the solution of an ordinary differential equation of the form

(3.52)
$$\frac{d\eta}{dt} = \omega(\eta) L[A'(\eta) - \frac{(A'(\eta), L'(\eta))}{\|\Psi'(\eta)\|^2} \Psi'(\eta)]$$

$$\eta(0,z) = z \in E.$$

The scalar function ω is Lipschitz continuous, $0 \le \omega(z) \le 1$, $\omega(z) = 1$ if $z \in M$ and A(z) is near c. Note that the argument of Proposition 3.23 shows $\Psi^{-1}(s)$ is a manifold for each s near 1, e.g. $|s-1| \le s_0$. The function $\omega(z) = 0$ if $|\Psi(z) - 1| > s_0$. Lastly $\omega(T_0z) = \omega(z)$ for all $\theta \in [0,2\pi]$.

Since the right hand side of (3.52) is Lipschitz continuous and is bounded by 1 (see [14] or [15]), there exists a solution of (3.52) defined for all teR and zeE.

Moreover $in(t,z) - zi \le 1$ for te[0,1], i.e. 3° holds. The form of (3.52) implies that n(t,M) = M for all teR. The properties of ω show that $n(1,\cdot)$ satisfies 1°-2°. Proposition 3.27 and a standard argument - see [13]-[15] imply 4°. To prove 5°, note that $p^+lA^+(z) = z^+$. Therefore integrating (3.52) yields:

(3.53)
$$P^{+}n(t,z) = (\exp \int_{0}^{t} \omega(n(s,z))ds)z^{+} - 0$$

$$-\int_{0}^{t} (\exp \int_{0}^{t} \omega(n(s,z))ds)\omega(n(\tau,z))(A^{+}(n(\tau,z),\Psi^{+}(n(\tau,z)))) ds$$

$$= 1\Psi^{+}(n(\tau,z))|_{0}^{-2}P^{+}L\Psi^{+}(n(\tau,z))d\tau.$$

Thus $P^{\dagger}n$ has the form stated in 5° . The compactness of P_2B^{\dagger} follows via Proposition 3.48 and an argument from [16] since P_2 and P^{\dagger} commute.

Remark 3.54: Let $A_g = \{z \in M \mid A(z) \le s\}$. If we replace $\omega(z)$ by $-\omega(z)$ in (3.52), the assertions of Proposition 3.51 still hold with 4^O replace by $\eta(1,A_{C+\epsilon}) \subseteq A_{C-\epsilon}$ and 5^O by $P^-\eta(1,z) = \beta^-(z)z^- + B^-(z)$ where $\beta^- \in C(M,[e^{-1},1])$ and P_2B^- is compact.

§4. Existence of a Solution

The proof of Theorem 1 will be completed in this section. The solution will be obtained as a critical point of $A|_{M}$ by a minimax argument. Then a simple regularity argument shows it is a classical solution of (HS). The following two lemmas pave the way for the definition of the critical value c.

Lemma 4.1: Let $M^+ \equiv M \cap E^+$ and set

$$\underline{\underline{\alpha}} = \inf_{z \in M} A(z)$$
.

Then $\underline{\alpha} > 0$.

<u>Proof:</u> By (2.2) for $z = z^+ \in E^+$, $A(z) = \frac{1}{2} \|z^+\|^2$. Since by Proposition 3.25 M is the boundary of a neighborhood of 0 in E, there is an r > 0 such that $\|z\| \le r$ implies z is interior to M. In particular for $z \in \partial B_r(0) \cap E^+$, $A(z) > \frac{1}{2} r^2$. Hence $\frac{a}{2} > \frac{1}{2} r^2$.

Next let L^+ be a two dimensional invariant subspace of E^+ . We further require that L^+ be such that there is a constant $a_1>0$ satisfying

for all $z \in E^0 \oplus E^- \oplus L^+$. To find such an L^+ , let e_1, \dots, e_n denote the usual basis in \mathbb{R}^n . Then we can take

 $E^0 = span\{(e_j, 0), (0, e_k) | 1 \le j, k \le n\}$,

 $E^+ \equiv span\{((j + 1)sin jt e_k - (1 + \frac{1}{j})cos jt e_k),$

 $(k + 1)\cos \ell t e_m, (1 + \frac{1}{j})\sin \ell t e_m) | 1 \le k, m \le n, \text{ and } j, \ell \in \mathbb{N} \}$

 $E^- \equiv \operatorname{span}\{((j+1)\sin jt e_k,(1+\frac{1}{j})\cos jt e_k),$

((l + 1)cos lt e_m , - (1 + $\frac{1}{\ell}$)sin lt e_m) | 1 < k,m < n and j₇l C N} ,

and L^+ = span{(sin t e₁,-cos t e₁),(cos t e₁,sin t e₁)}.

It is easy to verify that (4.2) holds.

Lemma 4.3: If $M^- \equiv M \cap (E^- \oplus E^0 \oplus L^+)$ and $\bar{\alpha} \equiv \sup_{z \in M} A(z)$,

then $\bar{\alpha} < \infty$.

<u>Proof:</u> By Proposition 3.26, M is bounded in $L^2(S^1, R^{2n})$. Therefore there is an $M_1 > 0$ such that $\|z\|_{L^2} \le M_1$ for all $z \in M$. In particular for $z = z^- + z^0 + z^+ \in M^-$, by (4.2) we have

$$||z^{+}||_{L^{2}} \le a_{1}M_{1}.$$

Since L⁺ is finite dimensional, there is a constant $a_2 > 0$ such that $\|z^{\dagger}\| \le a_2 \|z^{\dagger}\|_{L^2}$ for all $z^{\dagger} \in L^{\dagger}$. Hence

(4.5)
$$A(z) = \frac{1}{2} (\|z^{\dagger}\|^2 - \|z^{\dagger}\|^2) \le \frac{1}{2} (a_1 a_2 M_1)^2$$

for $z \in M^-$ and $\bar{\alpha} < \frac{1}{2} (a_1 a_2 M_1)^2$.

Now the class of sets that will be used to find a critical point of $\left.\mathbf{A}\right|_{M}$ can be introduced. Let

 $\Gamma \equiv \{h \in C(M,M) \mid 1^{\circ} h \text{ is equivariant,} \}$

 $2^0 h(z) = z \text{ if } A(z) \notin [0, \bar{\alpha} + 1],$

30 h(z) maps bounded sets to bounded sets

 4^0 P⁺h(z) = $\beta(z)z^+$ + B(z) where $\beta \in C(M,[1,\beta_0])$, $\beta_0 = \beta_0(h) > 1$, and P₂B(z) is compact}.

A critical value c of $A|_{M}$ can be produced by taking:

(4.6)
$$c \equiv \sup_{h \in \Gamma} \inf_{z \in M^+} A(h(z)) .$$

To see this, note first that id $\in \Gamma$. Hence by Lemma 4.1, $c \ge \underline{\alpha} > 0$. To prove that $c < \infty$, the following intersection theorem which is of independent interest is required. Theorem 4.7: Let $h \in \Gamma$. Then $h(M^+) \cap M^- \ne \emptyset$.

Proof: We will use a finite dimensional approximation argument. Let E_m^+ and E_m^- be E_m^- dimensional invariant subspaces of E^+,E^- respectively such that if $E_m = E_m^- \oplus E^0 \oplus E_m^+$, $\overline{\bigcup_{m \in \mathbb{N}}} = E$. Such subspaces can be written down explicitly using the basis for E given following $(4\cdot 2)$. Let P_m denote the orthogonal projector of E onto E_m^- . Set $h_m = P_m h \oplus C(M^+ \cap E_m, E_m^-)$. Note that by properties $1^{\circ}-2^{\circ}$ of Γ , h_m is equivariant and $h_m(z) = z$ on $E^0 \cap M$. By Proposition 2.2 of [17] (where we take f to be the orthogonal projector of E^+ onto the orthogonal complement of L^+ in E^+ composed with h_m), there is a point $z_m \oplus M^+ \cap E_m$ such that $h_m(z_m) \oplus E^- \oplus E^0 \oplus L^+$. We claim (z_m) is a bounded sequence. Otherwise $\|z_m\| + \infty$ along a subsequence. But then since $z_m \oplus E^+$, $A(z_m) = \frac{1}{2} \|z_m\|^2 + \infty$. Hence by property 2° of Γ , $h_m(z_m) = z_m$ for large m. Therefore $z_m \oplus M^+ \cap (E^- \oplus E^0 \oplus L^+) = M \cap L^+$. Since M is bounded in L^2 and L^+ is finite dimensional, (z_m) must be bounded in E, a contradiction.

Thus (z_m) is a bounded sequence. By property 3^O of Γ , $(h_m(z_m))$ is also bounded. Property 4^O of Γ implies that

(4.8)
$$q_{m} = \beta(z_{m})^{-1}(P_{2}P^{+}h_{m}(z_{m}) - P_{2}B(z_{m}))$$

Where $z_m=(p_m,q_m)$. The boundedness of (z_m) and compactness of P_2B show the second term on the right hand side of (4.8) has a convergent subsequence. The boundedness of $h_m(z_m)$ and the fact that $p^+h_m(z_m)$ lies in L^+ which is finite dimensional implies the first term on the right hand side of (4.8) also has a convergent subsequence. It follows then from (4.8) that q_m has a convergent subsequence in $W^{1,2}(S^1,\mathbb{R}^n)$. Therefore the same is true for $p_m=pq_m$ in $L^2(S^1,\mathbb{R}^n)$. Consequently $z_m+z\in M^+$ and by the continuity of h, $h_m(z_m)+h(z)\in M^-$. The Theorem is proved.

Corollary 4.9: $c \leq \overline{a} < \infty$.

Proof: By Theorem 4.7, $h(M^+) \cap M^- \neq \emptyset$ for any $h \in \Gamma$. Therefore for each $h \in \Gamma$,

$$\inf_{z \in M} A(h(z)) \leq \sup_{w \in M} A(w) = \frac{1}{\alpha}$$

via Lemma 4.3.

Now we can prove

Theorem 4.10: c is a critical value of A ...

<u>Proof:</u> If not, we can invoke Proposition 3.51 with $\overline{\epsilon} = \frac{1}{2} \min(1,\underline{\alpha})$ obtaining $\eta(1,\cdot) \in C(M,M)$ and satisfying $1^{\circ}-5^{\circ}$ of Proposition 3.51. But $1^{\circ}-3^{\circ}$, 5° , and our choice of $\overline{\epsilon}$ imply that $\eta(1,\cdot) \in \Gamma$ as is $\eta(1,h)$ for any $h \in \Gamma$. By 4° of Proposition 3.51,

$$(4.11) \eta(1, \cdot) : A_{c-\varepsilon} + A_{c+\varepsilon}.$$

Choose $h \in \Gamma$ so that

(4.12)
$$\inf_{z \in M^+} A(h(z)) \ge c - \epsilon .$$

By (4.11),

(4.13)
$$\inf_{z \in M} A(\eta(1,h(z))) \ge c + \varepsilon.$$

But since $\eta(1,h) \in \Gamma$, (4.6) shows

(4.14)
$$\inf_{z \in M^+} A(\eta(1,h(z)) \le c$$
,

a contradiction. Thus c is a critical value of A|M.

Now finally we can complete the

<u>Proof of Theorem 1:</u> Since c is a critical value of $A|_{M}$, there is a $\lambda \in \mathbb{R}$ and $z \in M$ such that A(z) = c and $A'(z) - \lambda \Psi'(z) = 0$, i.e.

(4.15)
$$\int_{0}^{2\pi} [(p \cdot \dot{Q}) + (P \cdot \dot{q}) - \frac{\lambda}{2\pi} (\ddot{H}_{p}(z) \cdot P) + (\ddot{H}_{q}(z) \cdot Q)] dt = 0$$

for all (P,Q) \in E. Equation (4.15) expresses the fact that z is a weak solution of (2.3). The argument of (3.28)-(3.30) and our lower bound for c show $\lambda > 0$. A simple regularity argument - see the proof of Theorem 3.3 of [18] - shows $z \in C^1(S^1, \mathbb{R}^{2n})$, i.e. z is a classical solution of (2.3). Therefore $H(z(t)) \equiv constant$ so $\Psi(z) = 1$ implies that $z(t) \in \mathcal{D}$. Lastly since $\lambda \neq 0$, making the change of time scale $t + \lambda t$ shows z is a $2\pi\lambda$ periodic solution of (HS). The proof is complete.

§5. A Dual Approach

In this section another existence proof will be given for a critical value of $A|_{M}$. This approach is "dual" to the previous one in the spirit of [19]. The critical value obtained in this section may differ from that given by (4.6).

The new critical value will as in §4 be obtained as a minimax. Let

 $\Lambda = \{g \in C(M,M) | g \text{ satisfies properties } 1^{O}-3^{O} \text{ of } \Gamma$ and $4^{O} P^{-}g = \beta^{-}(z)z^{-} + B^{+}(z) \text{ where } \beta^{-} \in C(M,[\beta_{\frac{1}{4}},1]),$ $\beta_{\frac{1}{4}} = \beta_{\frac{1}{4}}(g) > 0, \text{ and } P_{2}B^{-} \text{ is compact}\}.$

As in §4, there is an intersection theorem associated with Λ .

Theorem 5.1: If $g \in \Lambda$, $g(M^-) \cap M^+ \neq \emptyset$.

<u>Proof:</u> Set $g_m = P_m g \in C(M^- \cap E_m, E_m)$ where E_m and P_m are as in the proof of Theorem 4.7. By properties $1^{\circ}-2^{\circ}$ of Λ , g_m is equivariant and $g_m(z) = z$ on E° . Hence Proposition 2 of [17] can again be invoked – this time with f being the orthogonal projector of $E_m^- \oplus E^0 \oplus L^+$ onto $E_m^- \oplus E^0$ composed with g_m – to obtain $z_m \in M^- \cap E_m$ such that $g_m(z_m) \in E_m^+$. By Proposition 3.26, (z_m) is a bounded sequence in $L^2(S^1,R^{2n})$. Therefore by (4.2), (z_m^+) is bounded in L^2 and therefore in E since $z_m^+ \in L^+$ which is finite dimensional. Since E^0 is L^2 orthogonal to $E^- \oplus E^+$ via the definition of these spaces, (z_m^0) is bounded in E. We claim (z_m^-) is also bounded in E. If not,

$$A(z_m) = \frac{1}{2} (||z_m^+||^2 - ||z_m^-||^2) + -\infty.$$

But then by property 2° of Λ , $g_m(z_m) = z_m$ for large m so $z_m \in M^- \cap E^+ = M \cap L^+$. This implies $z_m = z_m^+$ for large m, $z_m^- = 0$, and (z_m^-) is a bounded sequence.

Since (z_m) is a bounded sequence, it possesses a subsequence which converges weakly in E to z e E. By property 4^O of Λ ,

(5.2)
$$P_{2}P^{-}g_{m}(z_{m}) = 0 = \beta^{-}(z_{m})q_{m}^{-} + P_{m}P_{2}B^{-}(z_{m})$$

and $P_m P_2 B^-(z_m)$ has a convergent subsequence in $W^{1,2}(S^1,\mathbb{R}^n)$. Hence so does q_m^- .

Therefore $p_m^- = D_{qm}^-$ does also in $L^2(S^1, \mathbb{R}^n)$. Since $\mathbb{E}^0 \oplus L^+$ is finite dimensional, it follows that $z_m + z$ in \mathbb{E} and $z \in M^-$. Since q is continuous, $g_m(z_m) + g(z) \in M \cap \mathbb{E}^+ = M^+$.

Now define

(5.3)
$$\widetilde{c} = \inf \sup_{g \in \Lambda} A(w) .$$

Theorem 5.4: \tilde{c} is a critical value of $A|_{M}$ with $\underline{\alpha} \leq \tilde{c} \leq \alpha$.

Proof: Since id $e \wedge \tilde{c} \leq \alpha$. Moreover by Theorem 5.1, if $g \in \Lambda$, $g(M^{-}) \cap M^{+} \neq \emptyset$.

Therefore

$$\tilde{c} > \inf_{M^+} A \equiv \underline{\alpha}$$
.

Finally using Remark 3.54, the proof that \tilde{c} is a critical value of $A|_{M}$ follows the same lines as the proof of Theorem 4.10 and we will omit it.

§6. An a Priori Bound for the Period

Theorem 1 establishes the existence of a periodic solution of (HS) as a critical point of $A|_{M}$. In this section, in a somewhat more general setting, an a priori bound will be obtained for the period of any periodic solution of (HS) in terms of A(z) and various constants determined from (H_1) - (H_3) . Writing (HS) in the form (2.3), the period is $2\pi\lambda$; hence our a priori bound is for λ .

Theorem 6.1: Suppose H satisfies (H_1) - (H_3) and $z \in C^1(S^1, \mathbb{R}^{2n})$ is a solution of (2.3) with $\lambda \neq 0$. Then there are constants $\overline{a} > \underline{a} > 0$ independent of z such that

$$(6.2) \underline{a}|\lambda(z)| \leq |\lambda| \leq \overline{a}|\lambda(z)|.$$

<u>Proof:</u> Without loss of generality we can assume λ and A(z) are positive. Writing (2.3) as

(6.3)
$$\dot{p} = -\lambda H_{q}(z)$$

$$\dot{q} = \lambda H_{p}(z) ,$$

equation (6.4) implies

(6.5)
$$A(z) = \lambda \int_{0}^{2\pi} p \cdot H_{p}(z) dt.$$

Consequently

(6.6)
$$A(z) \leq 2\pi\lambda \quad \max_{(\xi,\eta)=\zeta \in \mathcal{D}} \xi \cdot H_{p}(\zeta)$$

and this gives the lower bound for λ in (6.2).

Next from (6.3)

$$-\lambda \int_{0}^{2\pi} |H_{q}|^{2} dt = \int_{0}^{2\pi} \dot{p} \cdot H_{q} dt = -\int_{0}^{2\pi} p \cdot (H_{qq} \dot{q} + H_{qp} \dot{p}) dt$$

$$= -\lambda \int_{0}^{2\pi} p \cdot (H_{qq} H_{p} - H_{qp} H_{q}) dt$$

or

(6.7)
$$0 = \lambda \int_{0}^{2\pi} [|H_{q}|^{2} + (p \cdot (H_{qp}H_{q} - H_{qq}H_{p}))] dt.$$

Adding b times (6.7) to (6.5) gives

(6.8)
$$A(z) = \lambda \int_{0}^{2\pi} [p \cdot H_{p} + b|H_{q}|^{2} + (bp \cdot (H_{qp}H_{q} - H_{qq}H_{p}))]dt.$$

By (H_2) , there is a $\gamma > 0$ such that

$$|H_{\alpha}(0,q)| > \gamma$$
 if $(0,a) \in \mathcal{D}$.

Therefore there is a $\sigma > 0$ such that

(6.9)
$$|H_q(p,q)| > \frac{\gamma}{2}$$
 if $(p,q) \in \mathcal{D}$ and $|p| < \sigma$.

Making σ still smaller if necessary, it can be assumed that

(6.10)
$$|p \cdot (H_{qp}H_q - H_{qq}H_p)| \le \frac{\gamma^2}{8}$$
 if $(p,q) \in \mathcal{D}$ and $|p| \le \sigma$.

Writing (6.8) as

$$\frac{A(z)}{\lambda} \equiv I_1 + I_2$$

where I₁ denotes the integral of the right hand side of (6.8) over

 $\{t \in [0,2\pi] \mid |p(t)| \le \sigma\}$ and I_2 denotes the complementary integral, lower bounds will

be obtained for I_1, I_2 . By (6.9)-(6.11), if

$$\ell \equiv \text{meas}\{t \in [0,2\pi] | |p(t)| \leq \sigma\}$$
,

then

(6.12)
$$r_1 > b(\frac{\gamma^2}{4} - \frac{\gamma^2}{8}) \ell = b \frac{\gamma^2}{8} \ell .$$

To estimate I2, let

$$M_1 = \max_{z \in \mathcal{O}} |p \cdot (H_{qp}H_q - H_{qq}H_p)|$$

and

$$\omega(\sigma) = \frac{1}{2M_1} \min_{\mathbf{z} \in \mathcal{D}, |\mathbf{p}| > \sigma} \mathbf{p} \cdot \mathbf{H}_{\mathbf{p}}(\mathbf{z}) .$$

Then

(6.13)
$$I_2 > (2\omega(\sigma) - h)M_1(2\pi - l).$$

Choosing $b = \omega(\sigma)$ and combining (6.11)-(6.13) yields

(6.14)
$$\frac{A(z)}{\lambda} > \omega(\sigma) \left(\frac{\gamma^2}{8} + M_1(2\pi - \ell)\right)$$

$$\geq 2\pi\omega(\sigma)\min(\frac{\gamma^2}{8},M_1) = \kappa(\sigma)$$
.

Thus the upper bound for λ in (6.2) holds with $\bar{a} = \kappa(\sigma)^{-1}$.

Remark 6.15: The constant <u>a</u> in (6.2) depends only on C^1 bounds for H on $\mathcal D$ while \overline{a} depends on C^2 bounds for H. To obtain an existence theorem for (HS) when H is merely in C^1 , a better estimate for \overline{a} is needed. Let $W(z) \in \mathbb R^n$ such that W is C^1 in a neighborhood of $\mathcal D$. Then as in (6.7) we get

(6.16)
$$0 = \lambda \int_{0}^{2\pi} [(w(z) \cdot H_{q}) + (p \cdot (w_{p}(z)H_{q} - w_{q}(z)H_{p}))]dt.$$

Suppose W satisfies

(6.17)
$$W(z) \cdot H_q(z) > \frac{\gamma^2}{2} \text{ if } z = (0,q) \in \mathcal{D}$$
.

Arguing as in (6.8)-(6.15) then yields

(6.18)
$$\frac{A(z)}{\lambda} > \omega(\sigma)(\frac{\gamma^2}{8} \ell + \overline{M}_1(2\pi - \ell)) \equiv \kappa(\sigma)$$

where

Therefore we get an upper bound for λ of the desired type.

The existence of a W as in (6.17) follows from a result of Palais [20]. If E is a real Banach space, $0 \subset E$, and $\Phi \in C^1(0,\mathbb{R})$, then we E is a pseudogradient vector for Φ at $z \in 0$ if

$$(6.19)$$
 (i) $\|w\| \le 2\|\Phi^*(z)\|$

(ii)
$$\langle \Phi^{+}(z), w \rangle_{E^{+}, E} > \|\Phi^{+}(z)\|^{2}$$
.

If $\Phi \in C^1(E,R)$, $\widetilde{E} = \{z \in E^{|\Phi^1(z)} \neq 0\}$, W(z) is locally Lipschitz continuous, and W(z) is a pseudogradient vector for all $z \in \widetilde{E}$, W(z) is called a pseudogradient vector field on \widetilde{E} . Palais has proved [20].

Lemma 6.20: If $\Phi \in C^1(F, \mathbb{R})$, there exists a pseudogradient vector field for W on \widetilde{E} .

Choosing $F \equiv \mathbb{R}^{2n}$ and using $(H_1) - (H_2)$, it is easy to verify there exists such a W in our setting. Moreover by using a smooth partition of unity in the proof of Lemma 6.20 - see e.g. Lemma 1.6 of [15] - it can be assumed that W is smooth. Thus the estimate (6.2) holds even when $F \in C^1$.

§7. A More Refined Existence Theorem

The goal of this section is to prove

Theorem 7.1: Suppose H satisfies

(H₁) He
$$c^{1}(\mathbb{R}^{2n},\mathbb{R})$$

and $(H_2)-(H_3)$. Then (HS) has a periodic solution on $\mathcal{D} \equiv H^{-1}(1)$.

<u>Proof:</u> Since H satisfies (H_1) , one can find a sequence of functions $H_m \in C^2(\mathbb{R}^{2n}, \mathbb{R})$ such that H_m satisfies (H_2) - (H_3) and H_m converges to H in the C^1 norm uniformly on compact subsets of \mathbb{R}^{2n} . By Theorem 1, the equation

$$\dot{z}_{m} = \lambda_{m} J H_{mz}(z_{m})$$

has a 2π periodic solution z_m lying on $\mathcal{D}_m \equiv H_m^{-1}(1)$ with λ_m satisfying (6.2) with constants \underline{a}_m , \overline{a}_m . Moreover

(7.3)
$$\underline{\alpha}_{m} = \inf_{z \in M_{m}} A(z) \leq A(z_{m}) \leq \sup_{z \in M_{m}} A(z) = \overline{\alpha}_{m}.$$

We claim there exist constants $\alpha^* > \alpha_* > 0$ such that

$$\alpha_{\star} \leq \underline{\alpha}_{m} \leq \overline{\alpha}_{m} \leq \alpha^{\star}$$

for all $m \in N$ and constants $a^* > a_* > 0$ such that

$$a_{*}A(z_{m}) \leq \lambda_{m} \leq a^{*}A(z_{m})$$

for all large m e M. Assuming (7.4)-(7.5) for the moment, it follows that the sequence (λ_m) is bounded away from 0 and ∞ . Since \mathcal{D}_m is near \mathcal{D} for all large m, the functions z_m are bounded in $L^\infty(S^1,R^{2n})$. Hence by (7.2), (z_m) is bounded in $C^1(S^1,R^{2n})$. The Arzela-Ascoli Theorem and (7.2) then imply that (λ_m,z_m) converges in $\mathbb{R}\times C^1(S^1,R^{2n})$ to (λ,z) satisfying

$$\dot{z} = \lambda J H_2(z)$$
.

To complete the proof of Theorem 7.1, (7.4)-(7.5) must be verified. For the latter inequalities, the constant \underline{a}_m is determined from (6.6) with H replaced by H_m . Since $H_m + H$ in C^1 uniformly in a neighborhood of D, an a_* which works for all large $m \in M$ can be determined. The same reasoning, together with the proofs of Theorem 6.1 and Remark 6.15 supply an a^* independent of m provided that there exists a W(z)

satisfying (6.17) with $H=H_{\mathfrak{m}}$ but γ and W independent of $\mathfrak{m}.$ Since a γ exists such that

(7.6)
$$|H_{q}(z)| > \gamma^{2} \text{ for } z = (0,q) \in \mathcal{D},$$

the convergence of H_m to H implies there is a $k \in \mathbb{N}$ such that

(7.7)
$$H_{mq}(z) \cdot H_{lq}(z) > \frac{\gamma^2}{2}$$

for $z = (0,q) \in \mathcal{D}_m$ and for all m,l > k. Therefore W can be taken to be $H_{kq}(z)$.

Lastly to check (6.4), first note that by the construction of H in §1, it can be assumed that there are constants r_1, r_2 such that $H_m(z) > r_1|z|^2 - r_2$ for all $z \in \mathbb{R}^{2n}$ (independently of m). Therefore $z \in M_m = \Psi_m^{-1}(1)$ implies that $\|z\|_{L^2} < 2\pi (1+r_2)r_1^{-1}$. The proof of Lemma 4.3 then shows how to obtain α^* . To get α_* , we argue indirectly. If there were no such constant, then for each meN, there is a $\zeta_m \in M_m^+$ and such that $\|\zeta_m\| + 0$ as $m + \infty$. Suppose $\zeta_m = (\xi_m, \eta_m)$. Then $\eta_m + 0$ in $W^{1,2}(S^1, \mathbb{R}^n)$ and a fortior $\eta_m + 0$ in $L^\infty(S^1, \mathbb{R}^n)$ while $\xi_m + 0$ in $L^2(S^1, \mathbb{R}^n)$. Since $V_m(q) = U_m(q) = \rho^{-1}|q|^2$ for small q independently of m via the definition of H in §1, $m(\eta_m) + 0$ as $m + \infty$. By (K_4) of Proposition 1.3,

$$K_{m}(z) \le r(1 + |z|)|p|$$

where r is independent of m. Consequently $K_{\rm m}(\zeta_{\rm m}) + 0$ as m + ∞ . But then $1 = \Psi_{\rm m}(\zeta_{\rm m}) + 0$ as m + ∞ , a contradiction. The proof is complete.

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7. DISTRIBUTION STATEMENT (of the abstract entered in Block 20, if diffe	rent from Report)
8. SUPPLEMENTARY NOTES	
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